Nonexistence of unbiased test for high-dimensional linear model with fixed design and a Bayesian-motivated test

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Thursday 4th April, 2019

Abstract

This paper considers testing regression coefficients in high-dimensional linear model with fixed design matrix. We prove that there does not exist nontrivial unbiased test for this problem. Moreover, no test can guarantee nontrivial power even when the true model is largely deviate from the null hypothesis. Nevertheless, Bayesian methods can still produce tests with good average power behavior. We propose a new test statistic which is the limit of Bayes factors under normal distribution. The null distribution of the proposed test statistic is approximated by Lindeberg's replacement trick. Under certain conditions, the global asymptotic power function of the proposed test is also given. The finite sample performance of the proposed test is demonstrated via simulation studies.

Key words: High-dimensional test, Lindeberg method, Linear model, Unbiasedness.

1 Introduction

Consider linear regression model of the form

$$\mathbf{y} = \mathbf{X}_a \boldsymbol{\beta}_a + \mathbf{X}_b \boldsymbol{\beta}_b + \boldsymbol{\epsilon},\tag{1}$$

where $\mathbf{y} \in \mathbb{R}^n$ is the response, \mathbf{X}_a , \mathbf{X}_b are $n \times q$ and $n \times p$ design matrices, respectively, $\boldsymbol{\beta}_a \in \mathbb{R}^q$, $\boldsymbol{\beta}_b \in \mathbb{R}^p$ are unknown regression coefficients, and $\boldsymbol{\epsilon} \in \mathbb{R}^n$ is the vector of random errors. Here the

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design matrix \mathbf{X}_a contains the predictors that are known to have effects on the response, and our interest is to test if \mathbf{X}_b contains any useful predictors. That is, we would like to test the hypotheses

$$\mathcal{H}_0: \boldsymbol{\beta}_b = 0, \quad \text{v.s.} \quad \mathcal{H}_1: \boldsymbol{\beta}_b \neq 0.$$
 (2)

The conventional test for hypotheses (2) is the F-test which is also the likelihood ratio test under normal errors. However, the F-test is not well defined in high dimensional setting. In fact, if ϵ is normal distributed and Rank $[\mathbf{X}_a; \mathbf{X}_b] = n$, then the likelihood is unbounded under the alternative hypothesis. Thus, new test methodology is required in high-dimensional setting.

For the problem of testing hypotheses (2), two different high-dimensional settings have been extensively considered in the literature. One is the small p, large q setting. An important example of this setting is testing individual coefficients of a high-dimensional regression. See [3, 35, 27] for testing procedures in this setting. In this paper, however, we focus on the other setting, namely the large p, small q setting. In this case, there are just a few covariates, namely \mathbf{X}_a , are known to have effects on the response, while there remain a large number of covariates, namely \mathbf{X}_b , to be tested.

Many test procedures have been proposed in the large p, small q setting. Based on an empirical Bayes model, Goeman et al. [16] and Goeman et al. [17] proposed a score test statistic as well as a method to determine the critical value. This score test was further investigated by Lan et al. [26] and Lan et al. [25]. Based on U-statistics, Zhong and Chen [36] proposed a test for the case where $\mathbf{X}_a = \mathbf{1}_n$. Their method was then extended and improved by Wang and Cui [33] and Cui et al. [9], respectively. To accommodate outlying observations and heavy-tailed distributions, Feng et al. [14] proposed a rank-based test for the entire coefficients. Xu [34] modified the test of Feng et al. [14] and proposed a scalar invariant rank-based test. Janson et al. [20] utilized a convex optimization program to obtain an unbiased estimator of a signal to noise ratio with small variance, based on which a test is constructed. Apart from the afore mentioned tests, there is another line of research utilizing desparsified Lasso estimator; see [11] and the references therein.

Except for the test of Goeman et al. [16], most existing high dimensional tests adopted the random design assumption, that is, the rows of \mathbf{X}_b are considered as being generated from a super population. As noted by Lei et al. [29], assuming a fixed design or a random design could lead to qualitatively different inferential results and the former is preferable from a theoretical point of view. Hence we focus on the fixed design setting in this paper. Of course, our results are still valid for random design from a conditional inference perspective.

In Bayesian literature, many Bayesian tests have been proposed for hypothesis (2) in low-dimensional setting; see [21, 15, 37] and the references therein. So far, most existing Bayesian tests are not applicable in large p, small q setting. In principle, using Bayes factor as a frequentist test statistic in high-dimensional setting is a good strategy for at least two reasons. First, the Bayes factors corresponding to proper priors are always well defined, even if the likelihood is unbounded. Second, under mild conditions, tests based on Bayes factors are admissible (See, e.g., [28], Theorem

6.7.2). In fact, the test procedure of Goeman et al. [16] is motivated by Bayesian methods but is treated as a frequentist significance test.

The test statistic of Goeman et al. [16] is obtained by letting the prior magnitude tend to zero. Their goal is to achieve the highest average power under local alternative. In classical statistics, any reasonable test can detect large deviations from the null hypothesis and it is natural to consider only the local alternative. However, the situation is different for our problem. In fact, Theorem 1 shall show that no test can detect all large deviations from the null hypothesis in large p, small q setting. Consequently, it is not enough to consider only the local alternative. In this paper, we propose a new test statistic in large p, small q setting which is the limit of Bayes factors as the prior magnitude tend to infinity. The proposed test statistic is a ratio of quadratic forms of y. We give an approximation of the distribution of quadratic form using Lindeberg's replacement trick. Based on this approximated distribution, the critical value of the proposed test statistic is calculated by a one step iteration. We prove that the proposed test procedure is valid under weak conditions. In particular, we do not require that the test statistic is asymptotically normally distributed. Under certain conditions, we also derive the asymptotic power function of the proposed test and the test of Goeman et al. [16]. It is shown that the proposed test can detect the signals from more directions than the test of Goeman et al. [16]. Simulations are conducted to examine the performance of the proposed test.

The rest of the paper is organized as follows. In Section 2, we prove the nonexistence of unbiased test for hypotheses (2) in large p, small q setting. In Section 3, we propose a Bayesian-motivated test and study the theoretical properties of the proposed test procedure. The asymptotic power function is also derived. Section 4 contains the simulation results. Section 5 concludes the paper. The technical proofs are presented in Appendix.

2 Nonexistence of unbiased test

In this paper, we consider testing hypotheses (2) in large p, small q setting. More specifically, we make the following assumption throughout the paper.

Assumption 1. In the linear model (1), suppose $\boldsymbol{\epsilon} = (\epsilon_1, \dots, \epsilon_n)^{\top}$ where $\epsilon_1, \dots, \epsilon_n$ are independent indentically distributed (iid) with $E(\epsilon_1) = 0$ and $Var(\epsilon_1) = \phi^{-1}$. Furthermore, suppose q < n, $Rank(\mathbf{X}_a) = q$ and $Rank([\mathbf{X}_a; \mathbf{X}_b]) = n$.

Testing hypotheses (2) in large p, small q setting is a challenging problem. Goeman et al. [16] noticed that their test has negligible power for many alternatives and consequently is not an unbiased test. Biased tests are often regarded as problematic in classical statistics. However, the situation is very different for our problem. In fact, we shall show that under normal assumption, there is no nontrivial unbiased test in large p, small q setting.

Let $\mathbf{P}_a = \mathbf{X}_a (\mathbf{X}_a^{\top} \mathbf{X}_a)^{-1} \mathbf{X}_a^{\top}$ be the projection matrix onto the column space of \mathbf{X}_a and denote

$$\tilde{\mathbf{P}}_a = \mathbf{I}_n - \mathbf{P}_a. \text{ Define } \boldsymbol{\theta} = (\boldsymbol{\beta}_a^\top, \boldsymbol{\beta}_b^\top, \phi)^\top \text{ and}$$

$$\Theta_0 = \left\{ \boldsymbol{\theta} = (\boldsymbol{\beta}_a^\top, \boldsymbol{\beta}_b^\top, \phi)^\top : \boldsymbol{\beta}_a \in \mathbb{R}^q, \boldsymbol{\beta}_b = 0, \phi > 0 \right\},$$

$$\Theta = \left\{ \boldsymbol{\theta} = (\boldsymbol{\beta}_a^\top, \boldsymbol{\beta}_b^\top, \phi)^\top : \boldsymbol{\beta}_a \in \mathbb{R}^q, \boldsymbol{\beta}_b \in \mathbb{R}^p, \phi > 0 \right\}.$$

A test function $\varphi(\mathbf{y})$ of level α is a Borel measurable function from \mathbb{R}^n to \mathbb{R} such that $0 \leq \varphi(\mathbf{y}) \leq 1$ and $\mathrm{E}_{\boldsymbol{\theta}}[\varphi(\mathbf{y})] \leq \alpha$ for $\boldsymbol{\theta} \in \Theta_0$. Let $\lambda(\cdot)$ be the Lebesgue measure on \mathbb{R}^n . We have the following theorem.

Theorem 1. Suppose Assumption 1 holds. Furthermore, assume $\boldsymbol{\epsilon} \sim \mathcal{N}_n(0, \phi^{-1}\mathbf{I}_n)$. Let $\varphi(\mathbf{y})$ be a test function of level α , $0 < \alpha < 1$. If $\lambda(\{\mathbf{y} : \varphi(\mathbf{y}) \neq \alpha\}) > 0$, then for every $M \geq 0$, there exists a parameter $\boldsymbol{\theta} = (\boldsymbol{\beta}_a^\top, \boldsymbol{\beta}_b^\top, \phi)^\top \in \Theta$ such that $\phi \boldsymbol{\beta}_b^\top \mathbf{X}_b^\top \tilde{\mathbf{P}}_a \mathbf{X}_b \boldsymbol{\beta}_b > M$ and $\mathbf{E}_{\boldsymbol{\theta}}[\varphi(\mathbf{y})] < \alpha$.

It follows from Theorem 1 that there is no nontrivial unbiased test in large p, small q setting. But Theorem 1 provides more information. In Theorem 1, the quantity $\phi \boldsymbol{\beta}_b^{\top} \mathbf{X}_b^{\top} \tilde{\mathbf{P}}_a \mathbf{X}_b \boldsymbol{\beta}_b$ can be regarded as a signal to noise ratio. In any conventional sense, it is expected that the power of a reasonable test should tend to 1 as $\phi \boldsymbol{\beta}_b^{\top} \mathbf{X}_b^{\top} \tilde{\mathbf{P}}_a \mathbf{X}_b \boldsymbol{\beta}_b \to \infty$. Indeed, such a test can be easily constructed in low dimensional setting. However, Theorem 1 claims that it is impossible to find such a test in large p, small q setting.

Thus, testing hypotheses (2) in large p, small q setting is a difficult problem. Classical strategies for constructing tests may not be suitable since there is no nontrivial unbiased test nor test with guaranteed power. This motivates us to consider tests with good average power behavior. Then Bayesian methods are natural choices.

3 A Bayesian-motivated test

3.1 Test statistic

Bayes factor (see, e.g., [24]) is a commonly used tool for model comparison within the Bayesian framework. In our problem, suppose $\epsilon \sim \mathcal{N}_n(0, \phi^{-1}\mathbf{I}_n)$, the Bayes factor for hypotheses (2) is

$$B_{10} = \frac{\int d\mathcal{N}_n(\mathbf{X}_a\boldsymbol{\beta}_a + \mathbf{X}_b\boldsymbol{\beta}_b, \phi^{-1}\mathbf{I}_n)(\mathbf{y})\pi_1(\boldsymbol{\beta}_b, \boldsymbol{\beta}_a, \phi) \,\mathrm{d}\boldsymbol{\beta}_b \mathrm{d}\boldsymbol{\beta}_a \mathrm{d}\phi}{\int d\mathcal{N}_n(\mathbf{X}_a\boldsymbol{\beta}_a, \phi^{-1}\mathbf{I}_n)(\mathbf{y})\pi_0(\boldsymbol{\beta}_a, \phi) \,\mathrm{d}\boldsymbol{\beta}_a \mathrm{d}\phi},$$

where $d\mathcal{N}_n(\mu, \mathbf{\Sigma})(\mathbf{y})$ is the density function of a $\mathcal{N}_n(\mu, \mathbf{\Sigma})$ random vector evaluated at \mathbf{y} , $\pi_0(\boldsymbol{\beta}_a, \phi)$ and $\pi_1(\boldsymbol{\beta}_b, \boldsymbol{\beta}_a, \phi)$ are the prior densities under the null and alternative hypotheses, respectively. If B_{10} is large, the alternative hypothesis is preferred. It is known that if the true parameters indeed come from the specified prior distribution, then the Bayes factor is the likelihood ratio statistic and consequently, the test based on the Bayes factor is the most powerful test. We notice that in the frequentist framework, the Bayes factor can be treated as a test statistic, provided the critical value can be determined to preserve the test level.

The behavior of a Bayes factor largely depends on the choice of priors. In Bayesian literature, many priors have been considered for testing the coefficients of linear model. Popular priors include g-priors [30] and intrinsic priors [4]. Unfortunately, these priors are not well defined in large p, small q setting. Note that under the null hypothesis \mathcal{H}_0 , the model is low-dimensional. This allows us to impose the reference prior $\pi_0(\boldsymbol{\beta}_a, \phi) = c/\phi$, where c is a constant. Under \mathcal{H}_1 , write $\pi_1(\boldsymbol{\beta}_b, \boldsymbol{\beta}_a, \phi) = \pi_1(\boldsymbol{\beta}_b | \boldsymbol{\beta}_a, \phi) \pi_1(\boldsymbol{\beta}_a, \phi)$. For parameters $\boldsymbol{\beta}_a$ and ϕ , we consider the same prior as in \mathcal{H}_0 , that is $\pi_1(\boldsymbol{\beta}_a, \phi) = \pi_0(\boldsymbol{\beta}_a, \phi)$. For parameter $\boldsymbol{\beta}_b$, however, imposing the improper reference prior would not produce valid marginal density of \mathbf{y} . To make the marginal density of \mathbf{y} well defined, we consider the simple normal prior $p_1(\boldsymbol{\beta}_b | \boldsymbol{\beta}_a, \phi) = d\mathcal{N}_p(0, \kappa^{-1}\phi^{-1}\mathbf{I}_p)(\boldsymbol{\beta}_b)$, where $\kappa > 0$ is a hyperparameter. That is, we put the following priors,

$$\pi_0(\boldsymbol{\beta}_a, \phi) = \frac{c}{\phi}, \quad \pi_1(\boldsymbol{\beta}_b, \boldsymbol{\beta}_a, \phi) = \frac{c}{\phi} d\mathcal{N}_p \left(0, \frac{1}{\kappa \phi} \mathbf{I}_p \right) (\boldsymbol{\beta}_b). \tag{3}$$

Let $B_{10,\kappa}$ be the Bayes factor corresponding to the priors (3). It is straightforward to show that

$$2\log(B_{10,\kappa}) = p\log\kappa - \log|\mathbf{X}_b^{\top}\tilde{\mathbf{P}}_a\mathbf{X}_b + \kappa\mathbf{I}_p|$$
$$-(n-q)\log\left(1 - \frac{\mathbf{y}^{\top}\tilde{\mathbf{P}}_a\mathbf{X}_b\left(\mathbf{X}_b^{\top}\tilde{\mathbf{P}}_a\mathbf{X}_b + \kappa\mathbf{I}_p\right)^{-1}\mathbf{X}_b^{\top}\tilde{\mathbf{P}}_a\mathbf{y}}{\mathbf{y}^{\top}\tilde{\mathbf{P}}_a\mathbf{y}}\right).$$

Denote by $\tilde{\mathbf{P}}_a = \tilde{\mathbf{U}}_a \tilde{\mathbf{U}}_a^{\top}$ the rank decomposition of $\tilde{\mathbf{P}}_a$, where $\tilde{\mathbf{U}}_a$ is an $n \times (n-q)$ column orthogonal matrix. Let $\mathbf{X}_b^* = \tilde{\mathbf{U}}_a^{\top} \mathbf{X}_b$, $\mathbf{y}^* = \tilde{\mathbf{U}}_a^{\top} \mathbf{y}$. Let γ_i be the *i*th largest eigenvalue of $\mathbf{X}_b^* \mathbf{X}_b^{*\top}$, $i = 1, \ldots, n-q$. Denote by $\mathbf{X}_b^* = \mathbf{U}_b^* \mathbf{D}_b^* \mathbf{V}_b^{*\top}$ the singular value decomposition of \mathbf{X}_b^* , where \mathbf{U}_b^* , \mathbf{V}_b^* are $(n-q) \times (n-q)$ and $p \times (n-q)$ column orthogonal matrices, respectively, and $\mathbf{D}_b^* = \mathrm{diag}(\sqrt{\gamma_1}, \ldots, \sqrt{\gamma_{n-q}})$. Then we have

$$2\log(B_{10,\kappa}) = (n-q)\log\kappa - \sum_{i=1}^{n-q}\log(\gamma_i + \kappa) - (n-q)\log\left(1 - \frac{\mathbf{y}^{*\top}\mathbf{U}_b^*\mathbf{D}_b^* \left(\mathbf{D}_b^{*2} + \kappa\mathbf{I}_{n-q}\right)^{-1}\mathbf{D}_b^*\mathbf{U}_b^{*\top}\mathbf{y}^*}{\mathbf{y}^{*\top}\mathbf{y}^*}\right).$$

The main part of the above expression is

$$T_{\kappa} = \frac{\mathbf{y}^{*\top} \mathbf{U}_b^* \mathbf{D}_b^* \left(\mathbf{D}_b^{*2} + \kappa \mathbf{I}_{n-q} \right)^{-1} \mathbf{D}_b^* \mathbf{U}_b^{*\top} \mathbf{y}^*}{\mathbf{y}^{*\top} \mathbf{y}^*}.$$

A large value of T_{κ} supports the alternative hypothesis. Hence T_{κ} can be regarded as a frequentist test statistic. It remains to choose an appropriate hyperparameter κ . As Goeman et al. [16] noted, the priors should place most probability on the alternatives which are perceived as more interesting to detect. Their strategy is to let the prior magnitude tend to zero to obtain a test with good power behavior under local alternatives, that is, $\|\boldsymbol{\beta}_b\|$ is small. Note that the prior magnitude of $\boldsymbol{\beta}_b$ decreases with the increase of κ . In fact, if we let κ tends to infinity, the limit

$$\lim_{\kappa \to \infty} \kappa T_{\kappa} = \frac{\mathbf{y}^{*\top} \mathbf{X}_b^* \mathbf{X}_b^{*\top} \mathbf{y}^*}{\mathbf{y}^{*\top} \mathbf{y}^*}$$

is exactly the test statistic of Goeman et al. [16].

As implied by Theorem 1, however, testing hypotheses (2) in large p, small q setting is such a difficult problem that no test can guarantee reasonable power, even if the true model is largely deviate from the null hypothesis. In this case, the global power behavior may worth more consideration than the local power behavior. Thus, contrary to the strategy of Goeman et al. [16], we let κ tend to 0 to obtain a test with good average power behavior for large $\|\beta_b\|$. Note that the statistic T_{κ} degenerates to 1 as $\kappa \to 0$. Nevertheless, the scaled statistic $(T_{\kappa} - 1)/\kappa$ has a proper limit as $\kappa \to 0$. Thus, we proposed the following test statistic

$$T = \lim_{\kappa \to 0} \frac{T_{\kappa} - 1}{\kappa} = -\frac{\mathbf{y}^{*\top} (\mathbf{X}_b^* \mathbf{X}_b^{*\top})^{-1} \mathbf{y}^*}{\mathbf{y}^{*\top} \mathbf{y}^*}.$$

The null hypothesis will be rejected if T is large.

3.2 Critical value

To formulate a valid frequentist test, we need to determine the critical value of T. If ϵ were indeed normally distributed, then under the null hypothesis, $T \sim -(\sum_{i=1}^{n-q} \gamma_i^{-1} z_i^2)/(\sum_{i=1}^{n-q} z_i^2)$, where z_1, \ldots, z_{n-q} are iid $\mathcal{N}(0,1)$ random variables. In this case, the exact critical value can be easily obtained. In practice, however, normal distribution often serves as an approximation rather than the true distribution of data. We would like to derive an asymptotic valid critical value for T_n for general distributions of ϵ .

Under the null hypothesis,

$$T = -\frac{(\sqrt{\phi}\boldsymbol{\epsilon})^{\top} \tilde{\mathbf{U}}_{a} (\mathbf{X}_{b}^{*} \mathbf{X}_{b}^{*\top})^{-1} \tilde{\mathbf{U}}_{a}^{\top} (\sqrt{\phi}\boldsymbol{\epsilon})}{(\sqrt{\phi}\boldsymbol{\epsilon})^{\top} \tilde{\mathbf{U}}_{a} \tilde{\mathbf{U}}_{a}^{\top} (\sqrt{\phi}\boldsymbol{\epsilon})}.$$

The numerator and the denominator of the above expression are both quadratic forms of iid standardized random variables. Hence the key step towards the goal is to approximate the distribution of quadratic forms. The asymptotics of quadratic forms have been extensively studied; see, e.g., [1, 2, 12, 18, 22, 23]. Most existing work focused on the asymptotic normality of quadratic forms under certain conditions. However, normal distribution is just one of the possible limit distributions of quadratic forms. See Sevast'yanov [32] for a full characterization of the limit distributions of quadratic forms of normal random variables. We would like to approximate the distribution of quadratic forms under general conditions, not confined to the asymptotic normality case.

A general method for distribution approximation is Lindeberg's replacement trick, the technique used in the Lindeberg's original proof of central limit theorem. The idea of Lindeberg's replacement trick is to subsequently replace the random variables by suitable normal random variables; see, e.g., [31], Section III.4. Recently, this technique shows great power for many difficult problems; see [5] and the references therein. To approximate the distribution of quadratic forms, we shall apply Lindeberg's replacement trick to the squared terms and the cross-products terms, respectivey. Let $\mathscr{C}^4(\mathbb{R})$ denote the class of all bounded real functions on \mathbb{R} having bounded, continuous kth

derivatives, $0 \le k \le 4$. It is known that if $E f(Z_n) \to E f(Z)$ for every $f \in \mathscr{C}^4(\mathbb{R})$ then $Z_n \leadsto Z$; see, e.g., [31], Theorem 12 of Chapter III. We have the following approximation theorem.

Theorem 2. Let $\boldsymbol{\xi} = (\xi_1, \dots, \xi_n)^{\top}$, where ξ_i 's are iid random variable with $E \xi_1 = 0$, $Var(\xi_1) = 1$. Furthermore, suppose the distribution of ξ_1 is symmetric about 0 and has finite eighth moments. Let \mathbf{A} be an $n \times n$ symmetric matrix with elements $a_{i,j}$. Define

$$S = \frac{\boldsymbol{\xi}^{\top} \mathbf{A} \boldsymbol{\xi} - \operatorname{tr}(\mathbf{A})}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + (\operatorname{E}(\xi_1^4) - 3) \operatorname{tr}(\mathbf{A}^{\circ 2})}},$$

where $\mathbf{A}^{\circ 2} = \mathbf{A} \circ \mathbf{A}$ and \circ means Hadamard product. Let z_1, \ldots, z_n be iid random variables with distribution $\mathcal{N}(0,1)$ and $\check{z}_1, \ldots, \check{z}_n$ be iid random variables with distribution $\mathcal{N}(0,1)$ which are independent of ξ_1, \ldots, ξ_n . Let τ be a real number. Define

$$S_{\tau}^* = \frac{\tau \sum_{i=1}^n a_{i,i} \check{z}_i + 2 \sum_{1 \le i < j \le n} a_{i,j} z_i z_j}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + (\operatorname{E}(\xi_1^4) - 3) \operatorname{tr}(\mathbf{A}^{\circ 2})}}.$$

Then for every $f \in \mathscr{C}^4(\mathbb{R})$,

$$\begin{aligned}
&|\mathbf{E}f(S) - \mathbf{E}f(S_{\tau}^{*})| \\
&\leq \frac{|\mathbf{E}(\xi_{1}^{4}) - 1 - \tau^{2}|}{2} ||f^{(2)}||_{\infty} \frac{\operatorname{tr}(\mathbf{A}^{\circ 2})}{2\operatorname{tr}(\mathbf{A}^{2}) + (\mathbf{E}(\xi_{1}^{4}) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})} \\
&+ \frac{\max\left(\left|\mathbf{E}(\xi_{1}^{2} - 1)^{3}\right|, 12(\mathbf{E}(\xi_{1}^{4}) - 1)\right)}{6} ||f^{(3)}||_{\infty} \frac{\sum_{l=1}^{n} \left(|a_{l,l}| \sum_{i=1}^{n} a_{i,l}^{2}\right)}{\left(2\operatorname{tr}(\mathbf{A}^{2}) + (\mathbf{E}(\xi_{1}^{4}) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})\right)^{3/2}} \\
&+ \frac{16\operatorname{E}(\xi_{1}^{8}) + 80\operatorname{E}(\xi_{1}^{4}) + 3\tau^{4} + 96}{24} ||f^{(4)}||_{\infty} \frac{\sum_{l=1}^{n} \left(\sum_{i=1}^{n} a_{i,l}^{2}\right)^{2}}{\left(2\operatorname{tr}(\mathbf{A}^{2}) + (\mathbf{E}(\xi_{1}^{4}) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})\right)^{2}}.
\end{aligned} \tag{4}$$

Remark 1. If $\tau^2 = E(\xi_1^4) - 1$, the first term of the right hand side of (4) disappear. In practice, however, the quantity $E(\xi_1^4)$ is often unknown and τ^2 should be chosen as an estimator of $E(\xi_1^4) - 1$.

Remark 2. As noted in Chatterjee [6], Section 3.1, an almost necessary condition for the asymptotic normality of S is

$$\frac{\operatorname{tr}(\mathbf{A}^4)}{\left(2\operatorname{tr}(\mathbf{A}^2) + (\operatorname{E}(\xi_1^4) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})\right)^2} \to 0.$$
 (5)

On the other hand, it can be seen that the right hand side of (4) converges to 0 provided $\tau^2 = E(\xi_1^4) - 1$ and

$$\frac{\sum_{l=1}^{n} \left(\sum_{i=1}^{n} a_{i,l}^{2}\right)^{2}}{\left(2\operatorname{tr}(\mathbf{A}^{2}) + \left(\mathrm{E}(\xi_{1}^{4}) - 3\right)\operatorname{tr}(\mathbf{A}^{\circ 2})\right)^{2}} \to 0.$$
 (6)

It can be seen that (6) is much weaker than (5). For example, if $a_{i,j} = 1$, i = 1, ..., n, j = 1, ..., n and $E(\xi_1^4) = 3$, then the condition (6) holds but the condition (5) does not hold.

We now apply Theorem 2 to approximate the null distribution of the proposed statistic T. For $n \times n$ matrix \mathbf{A} and real number τ , let $F(x; \mathbf{A}, \tau)$ be the cumulative distribution function of $\tau \sum_{i=1}^{n} a_{i,i} \check{z}_i + 2 \sum_{1 \leq i < j \leq n} a_{i,j} z_i z_j$, where $z_1, \ldots, z_n, \check{z}_1, \ldots, \check{z}_n$ are iid $\mathcal{N}(0,1)$ random variables.

Under the null hypothesis, Theorem 2 implies that

$$\begin{aligned} & \Pr\left(T > x\right) \\ &= \Pr\left(\left(\sqrt{\phi}\boldsymbol{\epsilon}\right)^{\top} \left(-\tilde{\mathbf{U}}_{a}(\mathbf{X}_{b}^{*}\mathbf{X}_{b}^{*\top})^{-1}\tilde{\mathbf{U}}_{a}^{\top} - x\tilde{\mathbf{U}}_{a}\tilde{\mathbf{U}}_{a}^{\top}\right)\left(\sqrt{\phi}\boldsymbol{\epsilon}\right) > 0\right) \\ &\approx 1 - F\left(\operatorname{tr}\left((\mathbf{X}_{b}^{*}\mathbf{X}_{b}^{*\top})^{-1}\right) + (n - q)x; -\tilde{\mathbf{U}}_{a}(\mathbf{X}_{b}^{*}\mathbf{X}_{b}^{*\top})^{-1}\tilde{\mathbf{U}}_{a}^{\top} - x\tilde{\mathbf{U}}_{a}\tilde{\mathbf{U}}_{a}^{\top}, \sqrt{\phi^{2}\operatorname{E}(\epsilon_{1}^{4}) - 1}\right). \end{aligned}$$

Thus, in order to obtain a test with level α asymptotically, the ideal strategy is to find a consistent estimator $\hat{\tau}$ of $\sqrt{\phi^2 \operatorname{E}(\epsilon_1^4) - 1}$ and solve the critical value x from the equation

$$F\left(\operatorname{tr}(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^{-1} + (n-q)x; -\tilde{\mathbf{U}}_a(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^{-1}\tilde{\mathbf{U}}_a^\top - x\tilde{\mathbf{U}}_a\tilde{\mathbf{U}}_a^\top, \hat{\tau}\right) = 1 - \alpha.$$

However, solving this equation is not an easy task. For ease of implementation, we propose a one step iteration algorithm. We set the start point as

$$x^{(0)} = -\frac{\operatorname{tr}(\mathbf{X}_b^* \mathbf{X}_b^{*\top})^{-1}}{n - q},$$

which is chosen such that $\operatorname{tr}(-\tilde{\mathbf{U}}_a(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^{-1}\tilde{\mathbf{U}}_a^{\top}-x^{(0)}\tilde{\mathbf{U}}_a\tilde{\mathbf{U}}_a^{\top})=0$. Then the critical value obtained from the one step iteration is

$$x^{(1)} = \frac{F^{-1}\left(1 - \alpha; -\tilde{\mathbf{U}}_a(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^{-1}\tilde{\mathbf{U}}_a^{\top} - x^{(0)}\tilde{\mathbf{U}}_a\tilde{\mathbf{U}}_a^{\top}, \hat{\tau}\right) - \operatorname{tr}(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^{-1}}{n - q}.$$

As noted in Remark 1, here $\hat{\tau}^2$ should be a consistent estimator of $\phi^2 E(\epsilon_1^4) - 1$ under the null hypothesis.

Theorem 3. Suppose Assumption 1 holds. Furthermore, suppose the distribution of ϵ_1 is symmetric about 0 and has finite eighth moments. Suppose $E \epsilon_1^4 > \phi^{-2}$. Let

$$\mathbf{A} = -\tilde{\mathbf{U}}_a(\mathbf{X}_b^* \mathbf{X}_b^{*\top})^{-1} \tilde{\mathbf{U}}_a^{\top} - x^{(0)} \tilde{\mathbf{U}}_a \tilde{\mathbf{U}}_a^{\top} = -\tilde{\mathbf{U}}_a(\mathbf{X}_b^* \mathbf{X}_b^{*\top})^{-1} \tilde{\mathbf{U}}_a^{\top} + \frac{\operatorname{tr}\left((\mathbf{X}_b^* \mathbf{X}_b^{*\top})^{-1}\right)}{n-q} \tilde{\mathbf{U}}_a \tilde{\mathbf{U}}_a^{\top}.$$

Suppose as $n \to \infty$,

$$\frac{\max_{i,j} a_{i,j}^2}{\operatorname{tr}(\mathbf{A})^2} \to 0. \tag{7}$$

Let $\hat{\tau}^2$ be an consistent estimator of $\phi^2 E(\epsilon_1^4) - 1$ based on **X**, **y**. Then

$$\Pr\left(T > x^{(1)}\right) \to \alpha.$$

It remains to consistently estimate $\phi^2 \operatorname{E}(\epsilon_1^4) - 1$. Bai et al. [1] gave a consistent estimator of $\phi^2 \operatorname{E}(\epsilon_1^4) - 1$ based on the standardized residuals. Here we use a slightly different estimator which is based on the ordinary least squares residuals $\tilde{\boldsymbol{\epsilon}} = (\tilde{\epsilon}_1, \dots, \tilde{\epsilon}_n)^{\top} = \tilde{\mathbf{P}}_a \mathbf{y}$. From [1], Theorem 2.1,

$$E\left(\tilde{\boldsymbol{\epsilon}}^{\top}\tilde{\mathbf{P}}_{a}\tilde{\boldsymbol{\epsilon}}\right) = (n-q)\phi^{-1},$$

$$E\left(\sum_{i=1}^{n}\tilde{\epsilon}_{i}^{4}\right) = 3\phi^{-2}\operatorname{tr}(\tilde{\mathbf{P}}_{a}^{\circ 2}) + \left(E(\epsilon_{1}^{4}) - 3\phi^{-2}\right)\operatorname{tr}\left(\tilde{\mathbf{P}}_{a}^{\circ 2}\right)^{2}.$$

Then a moment estimator of $\phi^2 \operatorname{E}(\epsilon_1^4) - 1$ is

$$\hat{\tau}^2 = \frac{\frac{(n-q)^2 \sum_{i=1}^n \tilde{\epsilon}_i^4}{\left(\tilde{\boldsymbol{\epsilon}}^\top \tilde{\mathbf{P}}_a \tilde{\boldsymbol{\epsilon}}\right)^2} - 3\operatorname{tr}(\tilde{\mathbf{P}}_a^{\circ 2})}{\operatorname{tr}\left(\tilde{\mathbf{P}}_a^{\circ 2}\right)^2} + 2.$$

Proposition 1. Suppose the assumptions of Theorem 3 hold. Furthermore, suppose $q/n \to 0$. Then under the null hypothesis, $\hat{\tau}^2 \xrightarrow{P} \phi^2 E(\epsilon_1^4) - 1$.

We reject the null hypothesis if

$$T > x^{(1)}$$

This test procedure is asymptotically exact of size α under the conditions of Theorem 3 and Proposition 1.

3.3 Power analysis

In this section, we investigate the asymptotic power of the proposed test procedure as well as the test of Goeman et al. [16]. To make the expression of the asymptotic power functions tractable, we shall assume further conditions so that the test statistics are asymptotically normally distributed. Also, ϵ is assumed to be normally distributed so that we can obtain the global power function rather than only local power function.

To derive the asymptotic power of the proposed test and the test of Goeman et al. [16] simultaneously, we consider the general statistic $\mathbf{y}^{*\top}(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^k\mathbf{y}^*/\mathbf{y}^{*\top}\mathbf{y}^*$. In fact, the proposed test statistic corresponds to k=-1 while the test statistic of Goeman et al. [16] corresponds to k=1. Note that for any $x \in \mathbb{R}$,

$$\Pr\left(\frac{\mathbf{y}^{*\top}(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^k\mathbf{y}^*}{\mathbf{y}^{*\top}\mathbf{y}^*} \le x\right) = \Pr\left(\mathbf{y}^{*\top}\left((\mathbf{X}_b^*\mathbf{X}_b^{*\top})^k - x\mathbf{I}_{n-q}\right)\mathbf{y}^* \le 0\right).$$

Hence the asymptotic behavior of noncentral quadratic form will play a key role in our investigation. We have the following proposition.

Proposition 2. Let $Z = (z_1, \ldots, z_n)^{\top}$, where z_i 's are iid $\mathcal{N}(0,1)$ random variables. Let \mathbf{A} be an $n \times n$ symmetric matrix with elements $a_{i,j}$. Let $\mathbf{b} = (b_1, \ldots, b_n)^{\top}$ be an n dimensional vector.

If $\operatorname{tr}(\mathbf{A}^4)/\operatorname{tr}^2(\mathbf{A}^2) \to 0$, then

$$\frac{Z^{\top} \mathbf{A} Z + b^{\top} Z - \operatorname{tr}(\mathbf{A})}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + \|\mathbf{b}\|^2}} \rightsquigarrow \mathcal{N}(0, 1).$$

Remark 3. Proposition 2 does not impose any condition on **b**. This allows us to give the global asymptotic power function of tests. As the cost of this flexibility, we have to make the normal assumption.

Now we investigate the asymptotic behavior of $\mathbf{y}^{*\top}(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^k\mathbf{y}^*/\mathbf{y}^{*\top}\mathbf{y}^*$. Let $w_i = (\mathbf{V}_b^{*\top}\boldsymbol{\beta}_b)_i$ be the coordinate of β_b along the *i*th principal component direction of $\mathbf{X}_b^{*\top}\mathbf{X}_b^*$, $i=1,\ldots,n-q$. Let I be a random variable with uniform distribution on $\{1,\ldots,n-q\}$, that is, $\Pr(I=i)=i/(n-q)$, $i=1,\ldots,n-q$. It turns out that many quantities involved can be conveniently represented as the expectations with respect to I. For example, $\operatorname{tr}(\mathbf{X}_b^*\mathbf{X}_b^{*\top})^k/(n-q)=\sum_{i=1}^{n-q}\gamma_i^k/(n-q)=\mathrm{E}(\gamma_I^k)$.

Theorem 4. Suppose model (1) holds with $\epsilon \sim \mathcal{N}_n(0, \phi^{-1}\mathbf{I}_n)$. Let $k \neq 0$ be a fixed number. Suppose as $n \to \infty$, $n - q \to \infty$ and

$$\frac{\max_{1 \le i \le n-q} \left(\gamma_i^k - \mathcal{E}(\gamma_I^k)\right)^2}{(n-q)\operatorname{Var}(\gamma_I^k)} \to 0.$$
(8)

Then for any $x \in \mathbb{R}$,

$$\Pr\left(\frac{\mathbf{y}^{*\top} \left(\mathbf{X}_{b}^{*} \mathbf{X}_{b}^{*\top}\right)^{k} \mathbf{y}^{*}}{\mathbf{y}^{*\top} \mathbf{y}^{*}} \leq \mathrm{E}(\gamma_{I}^{k}) + \sqrt{\frac{2 \operatorname{Var}\left(\gamma_{I}^{k}\right)}{n - q}} x\right)$$

$$= \Phi\left(\frac{\left(\mathrm{E}(\gamma_{I} w_{I}^{2}) + \phi^{-1}\right) \sqrt{2 \operatorname{Var}\left(\gamma_{I}^{k}\right)} x - \sqrt{n - q} \operatorname{Cov}\left(\gamma_{I}^{k}, \gamma_{I} w_{I}^{2}\right)}{\sqrt{2 \phi^{-2} \operatorname{Var}(\gamma_{I}^{k}) + 4\phi^{-1} \operatorname{E}\left[\left(\gamma_{I}^{k} - \mathrm{E}(\gamma_{I}^{k}) - \sqrt{\frac{2 \operatorname{Var}(\gamma_{I}^{k})}{n - q}} x\right)^{2} \gamma_{I} w_{I}^{2}\right]}\right) + o(1),$$

where $\Phi(\cdot)$ is the cumulative distribution function of a standard normal random variable.

Under the conditions of Theorem 4, the proposed test should reject the null hypothesis when

$$\frac{\mathbf{y}^{*\top} \left(\mathbf{X}_b^* \mathbf{X}_b^{*\top}\right)^{-1} \mathbf{y}^*}{\mathbf{y}^{*\top} \mathbf{v}^*} \leq \mathrm{E}(\gamma_I^{-1}) + \sqrt{\frac{2 \operatorname{Var}\left(\gamma_I^{-1}\right)}{n-q}} \Phi^{-1}(\alpha),$$

and the asymptotic power function of the proposed test is

$$\Phi \left(\frac{\left(\mathrm{E}(\gamma_{I} w_{I}^{2}) + \phi^{-1} \right) \sqrt{2 \operatorname{Var} \left(\gamma_{I}^{-1} \right)} \Phi^{-1}(\alpha) + \sqrt{n - q} \operatorname{Cov} \left(-\gamma_{I}^{-1}, \gamma_{I} w_{I}^{2} \right)}{\sqrt{2 \phi^{-2} \operatorname{Var}(\gamma_{I}^{-1}) + 4 \phi^{-1} \operatorname{E} \left[\left(\gamma_{I}^{-1} - \mathrm{E}(\gamma_{I}^{-1}) - \sqrt{\frac{2 \operatorname{Var}(\gamma_{I}^{-1})}{n - q}} \Phi^{-1}(\alpha) \right)^{2} \gamma_{I} w_{I}^{2} \right]} \right). \tag{9}$$

On the other hand, the test of Goeman et al. [16] should reject the null hypothesis when

$$\frac{\mathbf{y}^{*\top}\mathbf{X}_{b}^{*}\mathbf{X}_{b}^{*\top}\mathbf{y}^{*}}{\mathbf{y}^{*\top}\mathbf{y}^{*}} > \mathrm{E}(\gamma_{I}) + \sqrt{\frac{2 \operatorname{Var}(\gamma_{I})}{n-q}} \Phi^{-1}(1-\alpha),$$

and the asymptotic power function of their test is

$$\Phi \left(\frac{\left(\operatorname{E}(\gamma_{I}w_{I}^{2}) + \phi^{-1} \right) \sqrt{2 \operatorname{Var}(\gamma_{I})} \Phi^{-1}(\alpha) + \sqrt{n - q} \operatorname{Cov}\left(\gamma_{I}, \gamma_{I}w_{I}^{2}\right)}{\sqrt{2 \phi^{-2} \operatorname{Var}(\gamma_{I}) + 4\phi^{-1} \operatorname{E}\left[\left(\gamma_{I} - \operatorname{E}(\gamma_{I}) + \sqrt{\frac{2 \operatorname{Var}(\gamma_{I})}{n - q}} \Phi^{-1}(\alpha) \right)^{2} \gamma_{I}w_{I}^{2} \right]} \right).$$
(10)

It can be seen from (9) and (10) that the power of the proposed test mainly depends on $Cov(-\gamma_I^{-1}, \gamma_I w_I^2)$ while the power of the test of Goeman et al. [16] mainly depends on $Cov(\gamma_I, \gamma_I w_I^2)$. Unfortunately, neither of these two quantities is positive definite. This fact is not surprising in view of Theorem 1. Nevertheless, $Cov(-\gamma_I^{-1}, \gamma_I w_I^2)$ and $Cov(\gamma_I, \gamma_I w_I^2)$ are positive definite if β_b is restricted in certain subspaces of \mathbb{R}^p .

Let d_1 be the maximum i such that $\gamma_i^{-1} < \mathrm{E}(\gamma_I^{-1})$. Let d_2 be the maximum i such that $\gamma_i > \mathrm{E}(\gamma_I)$. Then it can be seen that

$$Cov(-\gamma_I^{-1}, \gamma_I w_I^2) = \frac{1}{n-q} \sum_{i=1}^{n-q} (E(\gamma_I^{-1}) - \gamma_i^{-1}) \gamma_i w_i^2$$

is positive definite if $w_{d_1+1} = \cdots = w_{n-q} = 0$, while

$$Cov(\gamma_I, \gamma_I w_I^2) = \frac{1}{n-q} \sum_{i=1}^{n-q} (\gamma_i - \operatorname{E} \gamma_I) \gamma_i w_i^2$$

is positive definite if $w_{d_2+1} = \cdots = w_{n-q} = 0$. In other words, $Cov(-\gamma_I^{-1}, \gamma_I w_I^2)$ is positive definite if $\boldsymbol{\beta}_b$ belongs to the rank d_1 principal component subspace of $\mathbf{X}_b^{*\top} \mathbf{X}_b^*$, and $Cov(\gamma_I, \gamma_I w_I^2)$ is positive definite if $\boldsymbol{\beta}_b$ belongs to the rank d_2 principal component subspace of $\mathbf{X}_b^{*\top} \mathbf{X}_b^*$. Note that $\gamma_i > E(\gamma_I)$ implies that $\gamma_i^{-1} < (E(\gamma_I))^{-1} \le E(\gamma_I^{-1})$, where the last inequality follows from Jensen's inequality. It follows that $d_1 \ge d_2$. Hence the proposed test can always detect the signals from more directions than the test of Goeman et al. [16].

Now we consider a practical scenario where the proposed test is more favorable. Suppose \mathbf{X}_b is a random design matrix and the rows of \mathbf{X}_b are iid random vectors. Let $X_{b,i} \in \mathbb{R}^p$ be the transpose of the *i*th row of \mathbf{X}_b , $i=1,\ldots n$. Often in practice, the predictors are affected by a few common factors. Specifically, we assume $X_{b,i}$ has structre $X_{b,i} = \mathbf{B}\mathbf{f}_i + u_i$, $\mathbf{f}_i \in \mathbb{R}^K$ is the vector of common factors, K is factor number, \mathbf{B} is a $p \times K$ loading matrix, and $u_i \in \mathbb{R}^p$ is the vector of idiosyncratic components, $i=1,\ldots,n$. In this case, γ_1,\ldots,γ_K are typically much larger than $\gamma_{K+1},\ldots,\gamma_{n-q}$ (See, e.g., [13]). Note that $d_2 = K$ for large γ_1,\ldots,γ_K . In this case, the test of Goeman et al. [16] can only detect the signals from the rank K principal component subspace of $\mathbf{X}_b^{*\top}\mathbf{X}_b^*$. In some

situations, however, the effect of \mathbf{X}_b on \mathbf{y} is mainly made by the idiosyncratic components. In this case, $\boldsymbol{\beta}_b$ is almost orthogonal to the rank K principal component subspace of $\mathbf{X}_b^{*\top}\mathbf{X}_b^*$, and consequently, the test of Goeman et al. [16] will have poor behavior. On the other hand, the large eigenvalues $\gamma_1, \ldots, \gamma_K$ has little effect on d_1 . Hence the proposed test is expected to have good performance.

4 Numerical results

In this section, we conduct simulations to examine the empirical size and power of the proposed test (abbreviated as NEW) and compare it with the global test proposed by Goeman et al. [16] (abbreviated as GT) and the EigenPrism test prospoed by Janson et al. [20] (abbreviated as EP). Throughout our simulations, we take q = 10, p = 1000, $\alpha = 0.05$ and the empirical results are obtained from 3000 replications.

In our simulations, two distributions of ϵ_1 are considered, namely the chi-squared distribution $\epsilon_1 \sim (\chi^2(4)-4)/\sqrt{8}$ and the Student's t distribution $\epsilon_1 \sim t_9$. Define the signal to noise ratio (SNR) as SNR = $\sqrt{(n-q)\operatorname{Var}(\gamma_I)}\phi\|\beta_b\|^2/p$. We consider two structures of β_b : dense β_b and sparse β_b . In dense β_b setting, the coordinates of β_b are independently generated from the uniform distribution U(-c,c) where c is selected to reach certain SNR. In sparse β_b setting, we randomly generated from U(-c,c) where c is selected to reach certain SNR. The design matrices \mathbf{X}_a and \mathbf{X}_b are randomly generated beforehand and are fixed during each simulation. In all simulations, the elements of \mathbf{X}_a are iid from standard normal distribution. We consider three models of \mathbf{X}_b , as follows.

- Model I: the rows of \mathbf{X}_b are iid from $\mathcal{N}(0, \mathbf{\Sigma}_b)$, where $\mathbf{\Sigma}_b = \Gamma \Gamma^{\top}$ and Γ is a $p \times p$ matrix whose elements are iid generated from $\mathcal{N}(0, 1)$.
- Model II: the rows of \mathbf{X}_b are iid from $\mathcal{N}(0, \mathbf{\Sigma}_b)$, where $(\mathbf{\Sigma}_b)_{i,j} = 0.1$ if $i \neq j$ and $(\mathbf{\Sigma}_b)_{i,j} = 1$ otherwise.
- Model III: \mathbf{X}_b is generated as Model II but the observed design matrix is $\mathbf{F}\mathbf{B}^{\top} + \mathbf{X}_b$, where \mathbf{F} is an $n \times 2$ random matrix with iid $\mathcal{N}(0,1)$ entries and \mathbf{B} is a $p \times 2$ loading matrix with iid $\mathcal{N}(0,1)$ entries.

The simulation results are illustrated in Figures 1-3. From the results, it can be seen that all three tests can well maintain the significant level, although the EP test is a little conservative in some cases. Note that for n = 50, the powers of tests typically do not converge to 1 as SNR increases. This phenomenon, which is not surprising in view of Theorem 1, is eased when the sample size rises to 100. In all three simulations, the proposed test shows relatively good power behavior compared with its two competitors. Figure 3 shows that under Model III, the GT test has negligible power

while the proposed test has reasonable power. This verifies the argument made at the end of Section 3.3.

5 Conclusions

This paper is concerned with testing regression coefficients in large p, small q setting. It was proved that nontrivial unbiased test does not exists in this setting. We have proposed a Bayesian-motivated test statistic for high-dimensional linear model with fixed design. We proposed an approximation of the null distribution of the proposed test statistic which is then used to determine the critical value of the test statistic. Under weak conditions, we proved the proposed test procedure is asymptotically level α . Under certain conditions, we also derived the asymptotic power function of the proposed test.

The methodology of Goeman et al. [16] can also be applied in generalized linear models. However, it is not obvious how to generalize the proposed test and the test of Janson et al. [20] to generalized linear models. More fundamentally, it is unclear if results like Theorem 1 hold for generalized linear model.

In Theorem 2, we assumed that the distribution of ϵ_1 is symmetric about 0. This condition is also assumed by Bai et al. [1] in the study of central limit theorem of quadratic form. Without this condition, our proof of Theorem 2 is not valid. We think it will be an interesting and useful work to relax the symmetric condition in Theorem 2.

Acknowledgments

This work was supported by the National Natural Science Foundation of China under Grant No. 11471035.

Appendix

Lemma 1. Under the assumptions of Theorem 1, if there exists a Borel set $G \subset \mathbb{R}^n$ and a number M > 0 such that

$$\int_{\mathbb{R}^n} \varphi(\mathbf{y}) d\mathcal{N}_n(\mu, \phi^{-1} \mathbf{I}_n)(\mathbf{y}) \, \mathrm{d}\mathbf{y} \ge \alpha \quad \text{for all } \mu \in G \text{ and } \phi > M,$$

then $\varphi(\mathbf{y})\mathbf{1}_G(\mathbf{y}) \geq \alpha \mathbf{1}_G(\mathbf{y})$ a.e. λ .

Proof. We prove the claim by contradiction. Suppose $\lambda(\{\mathbf{y}:\varphi(\mathbf{y})<\alpha\}\cap G)>0$. Then there exists a sufficiently small $0<\eta<\alpha$, such that $\lambda(\{\mathbf{y}:\varphi(\mathbf{y})<\alpha-\eta\}\cap G)>0$. We denote $E:=\{\mathbf{y}:\varphi(\mathbf{y})<\alpha-\eta\}\cap G$. From Lebesgue density theorem [8, Corollary 6.2.6], there exists a

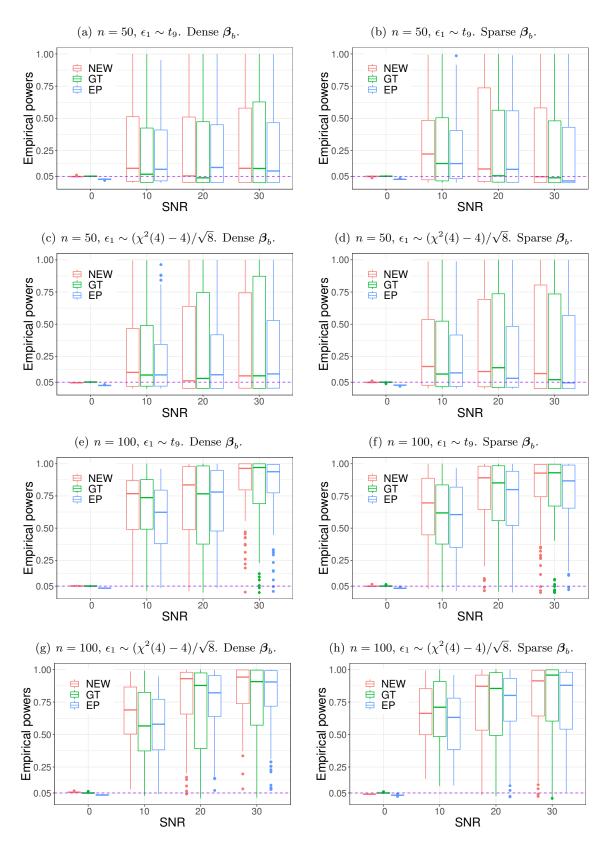


Figure 1: Box plots of the empirical powers based on 100 independently generated β_b . \mathbf{X}_b is generated by Model I.

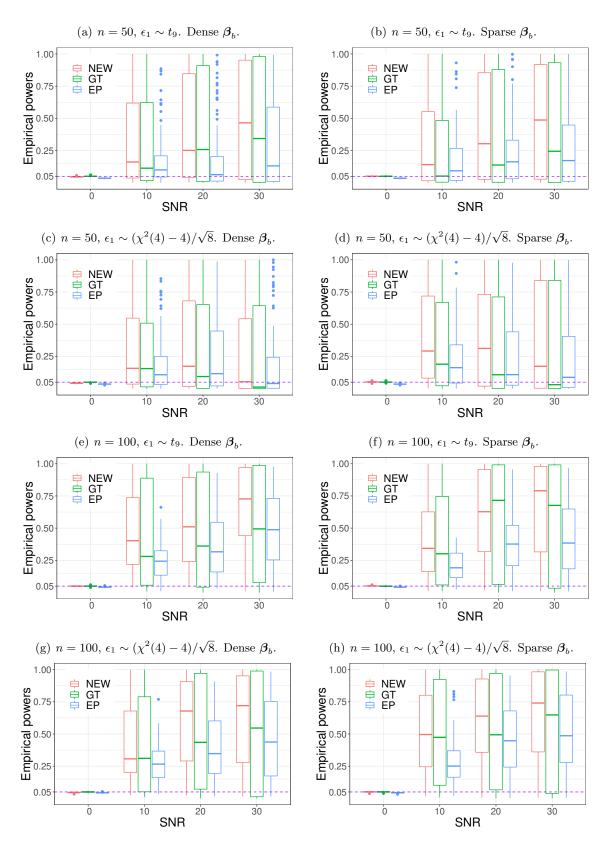


Figure 2: Box plots of the empirical powers based on 100 independently generated β_b . \mathbf{X}_b is generated by Model II.

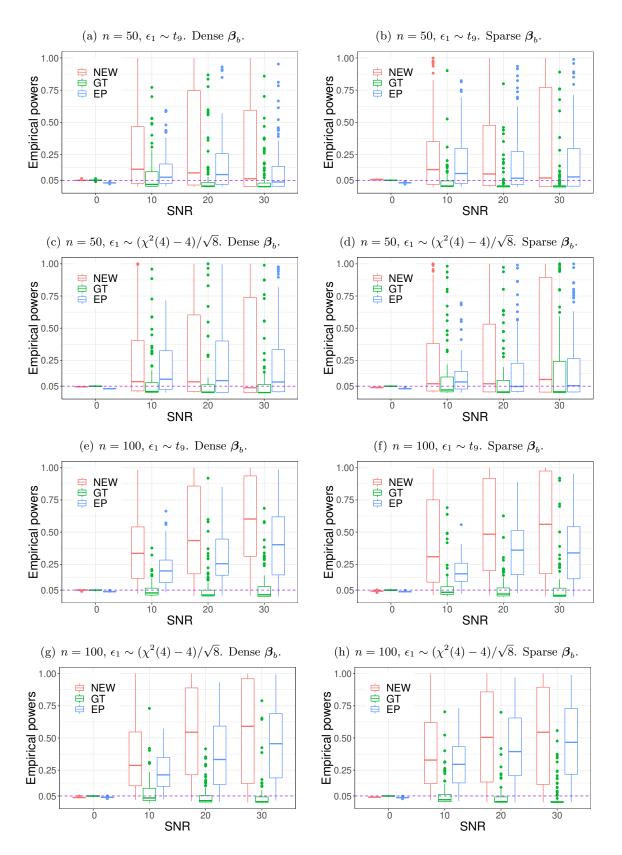


Figure 3: Box plots of the empirical powers based on 100 independently generated β_b . \mathbf{X}_b is generated by Model III.

point $z \in E$, such that, for any $\varepsilon > 0$ there is a $\delta_{\varepsilon} > 0$ such that for any $0 < \delta' < \delta_{\varepsilon}$,

$$\left| \frac{\lambda(E^{\complement} \cap C_{\delta'})}{\lambda(C_{\delta'})} \right| < \varepsilon,$$

where $C_{\delta'} = \prod_{i=1}^{n} [z_i - \delta', z_i + \delta']$. We put

$$\varepsilon = \left(\frac{\sqrt{\pi}}{\sqrt{2}\Phi^{-1}\left(1 - \frac{\eta}{6n}\right)}\right)^n \frac{\eta}{3}.$$

Then for any $\phi > M$ and $0 < \delta' < \delta_{\varepsilon}$,

$$\alpha \leq \int_{\mathbb{R}^{n}} \varphi(\mathbf{y}) d\mathcal{N}_{n}(z, \phi^{-1} \mathbf{I}_{n})(\mathbf{y}) \, d\mathbf{y}$$

$$= \int_{E \cap C_{\delta'}} \varphi(\mathbf{y}) d\mathcal{N}_{n}(z, \phi^{-1} \mathbf{I}_{n})(\mathbf{y}) \, d\mathbf{y} + \int_{E^{\complement} \cap C_{\delta'}} \varphi(\mathbf{y}) d\mathcal{N}_{n}(z, \phi^{-1} \mathbf{I}_{n})(\mathbf{y}) \, d\mathbf{y}$$

$$+ \int_{C_{\delta'}^{\complement}} \varphi(\mathbf{y}) d\mathcal{N}_{n}(z, \phi^{-1} \mathbf{I}_{n})(\mathbf{y}) \, d\mathbf{y}$$

$$\leq \alpha - \eta + \int_{E^{\complement} \cap C_{\delta'}} d\mathcal{N}_{n}(z, \phi^{-1} \mathbf{I}_{n})(\mathbf{y}) \, d\mathbf{y} + \int_{C_{\delta'}^{\complement}} d\mathcal{N}_{n}(z, \phi^{-1} \mathbf{I}_{n})(\mathbf{y}) \, d\mathbf{y}$$

$$\leq \alpha - \eta + \left(\frac{\phi}{2\pi}\right)^{n/2} \lambda(E^{\complement} \cap C_{\delta'}) + 2n \left(1 - \Phi(\sqrt{\phi}\delta')\right)$$

$$\leq \alpha - \eta + \left(\frac{\phi}{2\pi}\right)^{n/2} \varepsilon(2\delta')^{n} + 2n \left(1 - \Phi(\sqrt{\phi}\delta')\right)$$

$$= \alpha - \eta + \left(\frac{\sqrt{\phi}\delta'}{\Phi^{-1}\left(1 - \frac{\eta}{6n}\right)}\right)^{n} \frac{\eta}{3} + 2n \left(1 - \Phi(\sqrt{\phi}\delta')\right).$$

In the last inequality, we put δ' small enough such that

$$\left(\frac{\Phi^{-1}\left(1-\frac{\eta}{6n}\right)}{\delta'}\right)^2 > M,$$

and put

$$\phi = \left(\frac{\Phi^{-1}\left(1 - \frac{\eta}{6n}\right)}{\delta'}\right)^2.$$

Then we obtain the contradiction $\alpha \leq \alpha - (2/3)\eta$. This completes the proof.

Proof of Theorem 1. We prove the claim by contradiction. Suppose there exists an $M \geq 0$ such that

$$\int_{\mathbb{R}^n} \varphi(\mathbf{y}) d\mathcal{N}_n \left(\mathbf{X}_a \boldsymbol{\beta}_a + \mathbf{X}_b \boldsymbol{\beta}_b, \phi^{-1} \mathbf{I}_n \right) (\mathbf{y}) \, \mathrm{d}\mathbf{y} \ge \alpha$$

for every $\boldsymbol{\beta}_a \in \mathbb{R}^q$, $\boldsymbol{\beta}_b \in \mathbb{R}^p$, $\phi > 0$ satisfying $\phi \boldsymbol{\beta}_b^{\top} \mathbf{X}_b^{\top} \tilde{\mathbf{P}}_a \mathbf{X}_b \boldsymbol{\beta}_b > M$. Note that for any h > 0

$$\left\{(\boldsymbol{\beta}_b^\top, \phi)^\top: \boldsymbol{\beta}_b^\top \mathbf{X}_b^\top \tilde{\mathbf{P}}_a \mathbf{X}_b \boldsymbol{\beta}_b > h\sqrt{M}, \phi > h^{-1}\sqrt{M}\right\} \subset \left\{(\boldsymbol{\beta}_b^\top, \phi)^\top: \phi \boldsymbol{\beta}_b^\top \mathbf{X}_b^\top \tilde{\mathbf{P}}_a \mathbf{X}_b \boldsymbol{\beta}_b > M\right\}.$$

Then Lemma 1 implies that for any h > 0, $\varphi(\mathbf{y})\mathbf{1}_{G_h}(\mathbf{y}) \ge \alpha \mathbf{1}_{G_h}(\mathbf{y})$ a.e. λ , where

$$G_h = \left\{ \mathbf{X}_a \boldsymbol{\beta}_a + \mathbf{X}_b \boldsymbol{\beta}_b : \boldsymbol{\beta}_a \in \mathbb{R}^q, \boldsymbol{\beta}_b^{\top} \mathbf{X}_b^{\top} \tilde{\mathbf{P}}_a \mathbf{X}_b \boldsymbol{\beta}_b > h \sqrt{M}
ight\}.$$

It can be seen that $\lambda(\{\bigcup_{n=1}^{\infty} G_{1/n}\}^{\complement}) = 0$. Hence $\varphi(\mathbf{y}) \geq \alpha$ a.e. λ . On the other hand, since $\varphi(\mathbf{y})$ is a level α test, for every $\phi > 0$,

$$\int_{\mathbb{R}^n} [\varphi(\mathbf{y}) - \alpha] d\mathcal{N}_n(0, \phi^{-1} \mathbf{I}_n)(\mathbf{y}) \, d\mathbf{y} \le 0.$$
(11)

Note that the integrand of (11) is nonnegative. It follows that $\varphi(\mathbf{y}) = \alpha$ a.s. λ , a contradiction. This completes the proof.

Proof of Theorem 2. Let

$$\tilde{a}_{i,j} := \frac{a_{i,j}}{\sqrt{2\operatorname{tr}(\mathbf{A}^2) + (\operatorname{E}(\xi_1^4) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})}}.$$

Then

$$S = \sum_{i=1}^{n} \tilde{a}_{i,i}(\xi_i^2 - 1) + 2\sum_{1 \le i < j \le n} \tilde{a}_{i,j}\xi_i\xi_j, \quad S_{\tau}^* = \tau \sum_{i=1}^{n} \tilde{a}_{i,i}\check{z}_i + 2\sum_{1 \le i < j \le n} \tilde{a}_{i,j}z_iz_j.$$

For $l = 1, \ldots, n$, define

$$\begin{split} S_{l} &= \sum_{i=1}^{l-1} \tilde{a}_{i,i}(\xi_{i}^{2} - 1) + \tau \sum_{i=l+1}^{n} \tilde{a}_{i,i} \check{z}_{i} \\ &+ 2 \sum_{1 \leq i < j \leq l-1} \tilde{a}_{i,j} \xi_{i} \xi_{j} + 2 \sum_{i=1}^{l-1} \sum_{j=l+1}^{n} \tilde{a}_{i,j} \xi_{i} z_{j} + 2 \sum_{l+1 \leq i < j \leq n} \tilde{a}_{i,j} z_{i} z_{j}, \\ h_{l} &= \tilde{a}_{l,l}(\xi_{l}^{2} - 1) + 2 \sum_{i=1}^{l-1} \tilde{a}_{i,l} \xi_{i} \xi_{l} + 2 \sum_{i=l+1}^{n} \tilde{a}_{i,l} z_{i} \xi_{l}, \\ g_{l} &= \tau \tilde{a}_{l,l} \check{z}_{l} + 2 \sum_{i=1}^{l-1} \tilde{a}_{i,l} \xi_{i} z_{l} + 2 \sum_{i=l+1}^{n} \tilde{a}_{i,l} z_{i} z_{l}. \end{split}$$

It can be seen that for $l=2,\ldots,n$, $S_{l-1}+h_{l-1}=S_l+g_l$, and $S=S_n+h_n$, $S_1+g_1=S_{\tau}^*$. Thus, for any $f\in\mathscr{C}^4(\mathbb{R})$,

$$|E f(S) - E f(S_{\tau}^{*})|$$

$$= |E f(S_{n} + h_{n}) - E f(S_{1} + g_{1})|$$

$$= \left| \sum_{l=2}^{n} (E f(S_{l} + h_{l}) - E f(S_{l-1} + h_{l-1})) + E f(S_{1} + h_{1}) - E f(S_{1} + g_{1}) \right|$$

$$= \left| \sum_{l=1}^{n} E f(S_{l} + h_{l}) - E f(S_{l} + g_{l}) \right|.$$

Apply Taylor's theorem, for l = 1, ..., n,

$$f(S_l + h_l) = f(S_l) + \sum_{k=1}^{3} \frac{1}{k!} h_l^k f^{(k)}(S_l) + \frac{1}{24} h_l^4 f^{(4)}(S_l + \theta_1 h_l),$$

$$f(S_l + g_l) = f(S_l) + \sum_{k=1}^{3} \frac{1}{k!} g_l^k f^{(k)}(S_l) + \frac{1}{24} g_l^4 f^{(4)}(S_l + \theta_2 g_l),$$

where $\theta_1, \theta_2 \in [0, 1]$. Thus,

$$|\operatorname{E} f(S_l + h_l) - \operatorname{E} f(S_l + g_l)|$$

$$\leq \left| \sum_{k=1}^{3} \frac{1}{k!} \operatorname{E} f^{(k)}(S_l) \operatorname{E}_l(h_l^k - g_l^k) \right| + \frac{1}{24} ||f''''||_{\infty} \left(\operatorname{E}(h_l^4) + \operatorname{E}(g_l^4) \right),$$

where E_l denotes taking expectation with respect to ξ_l, z_l, \check{z}_l . It is straightforward to show that

$$\begin{split} & \mathbf{E}_l(h_l - g_l) = 0, \\ & \mathbf{E}_l(h_l^2 - g_l^2) = \left(\mathbf{E}(\xi_1^4) - 1 - \tau^2 \right) \tilde{a}_{l,l}^2, \\ & \mathbf{E}_l(h_l^3 - g_l^3) = \mathbf{E}(\xi_1^2 - 1)^3 \tilde{a}_{l,l}^3 + 12(\mathbf{E}(\xi_1^4) - 1) \tilde{a}_{l,l} \left(\sum_{i=1}^{l-1} \tilde{a}_{i,l} \xi_i + \sum_{i=l+1}^n \tilde{a}_{i,l} z_i \right)^2. \end{split}$$

Thus,

$$|\operatorname{E} f(S_{l} + h_{l}) - \operatorname{E} f(S_{l} + g_{l})|$$

$$\leq \frac{1}{2} \|f^{(2)}\|_{\infty} \left| \operatorname{E}(\xi_{1}^{4}) - 1 - \tau^{2} \right| \tilde{a}_{l,l}^{2}$$

$$+ \frac{1}{6} \|f^{(3)}\|_{\infty} \left(\left| \operatorname{E}(\xi_{1}^{2} - 1)^{3} \right| |\tilde{a}_{l,l}^{3}| + 12(\operatorname{E}(\xi_{1}^{4}) - 1) |\tilde{a}_{l,l}| \operatorname{E} \left(\sum_{i=1}^{l-1} \tilde{a}_{i,l} \xi_{i} + \sum_{i=l+1}^{n} \tilde{a}_{i,l} z_{i} \right)^{2} \right)$$

$$+ \frac{1}{24} \|f^{(4)}\|_{\infty} \left(\operatorname{E}(h_{l}^{4}) + \operatorname{E}(g_{l}^{4}) \right)$$

$$\leq \frac{\left| \operatorname{E}(\xi_{1}^{4}) - 1 - \tau^{2} \right|}{2} \|f^{(2)}\|_{\infty} \tilde{a}_{l,l}^{2} + \frac{\max\left(\left| \operatorname{E}(\xi_{1}^{2} - 1)^{3} \right|, 12(\operatorname{E}(\xi_{1}^{4}) - 1) \right)}{6} \|f^{(3)}\|_{\infty} |\tilde{a}_{l,l}| \sum_{i=1}^{n} \tilde{a}_{i,l}^{2}$$

$$+ \frac{1}{24} \|f^{(4)}\|_{\infty} \left(\operatorname{E}(h_{l}^{4}) + \operatorname{E}(g_{l}^{4}) \right) .$$

$$(12)$$

Now we bound $E(h_l^4)$ and $E(g_l^4)$. By direct calculation,

$$\begin{split} \mathbf{E}(h_{l}^{4}) &= \mathbf{E}(\xi_{1}^{2}-1)^{4} \tilde{a}_{l,l}^{4} + 24 \, \mathbf{E}[\xi_{1}^{2}(\xi_{1}^{2}-1)^{2}] \tilde{a}_{l,l}^{2} \, \mathbf{E}\left(\sum_{i=1}^{l-1} \tilde{a}_{i,l} \xi_{i} + \sum_{i=l+1}^{n} \tilde{a}_{i,l} z_{i}\right)^{2} \\ &+ 16 \, \mathbf{E}(\xi_{1}^{4}) \, \mathbf{E}\left(\sum_{i=1}^{l-1} \tilde{a}_{i,l} \xi_{i} + \sum_{i=l+1}^{n} \tilde{a}_{i,l} z_{i}\right)^{4} \\ &= \mathbf{E}(\xi_{1}^{2}-1)^{4} \tilde{a}_{l,l}^{4} + 24 \, \mathbf{E}[\xi_{1}^{2}(\xi_{1}^{2}-1)^{2}] \tilde{a}_{l,l}^{2} \left(\left(\sum_{i=1}^{n} \tilde{a}_{i,l}^{2}\right) - \tilde{a}_{l,l}^{2}\right) \\ &+ 16 \, \mathbf{E}(\xi_{1}^{4}) \left(\left(\mathbf{E}(\xi_{1}^{4}) - 3\right) \sum_{i=1}^{l-1} \tilde{a}_{i,l}^{4} + 3 \left(\left(\sum_{i=1}^{n} \tilde{a}_{i,l}^{2}\right) - \tilde{a}_{l,l}^{2}\right)^{2}\right). \end{split}$$

To upper bound the above quantity, we use the facts $24 \operatorname{E}[\xi_1^2(\xi_1^2-1)^2] \leq 2(16 \operatorname{E}(\xi_1^2-1)^4 + (9/4)\operatorname{E}(\xi_1^4))$, $\operatorname{E}(\xi_1^2-1)^4 \leq \operatorname{E}(\xi_1^8)$ and

$$\left(\mathbf{E}(\xi_1^4) - 3 \right) \sum_{i=1}^{l-1} \tilde{a}_{i,l}^4 \le \left(\mathbf{E}(\xi_1^4) - 1 \right) \sum_{i=1}^{l-1} \tilde{a}_{i,l}^4 \le \left(\mathbf{E}(\xi_1^4) - 1 \right) \left(\left(\sum_{i=1}^n \tilde{a}_{i,l}^2 \right) - \tilde{a}_{l,l}^2 \right)^2.$$

Then we obtain the bound

$$E(h_l^4) \le \left(16 E(\xi_1^8) + 32 E(\xi_1^4)\right) \left(\sum_{i=1}^n \tilde{a}_{i,l}^2\right)^2.$$
(13)

Similarly, we have

$$E(g_l^4) \le \left(48 E(\xi_1^4) + 3\tau^4 + 96\right) \left(\sum_{i=1}^n \tilde{a}_{i,l}^2\right)^2. \tag{14}$$

Combining (12), (13) and (14) yields

$$\sum_{l=1}^{n} |\operatorname{E} f(S_{l} + h_{l}) - \operatorname{E} f(S_{l} + g_{l})|$$

$$\leq \frac{\left|\operatorname{E}(\xi_{1}^{4}) - 1 - \tau^{2}\right|}{2} ||f^{(2)}||_{\infty} \sum_{l=1}^{n} \tilde{a}_{l,l}^{2}$$

$$+ \frac{\max\left(\left|\operatorname{E}(\xi_{1}^{2} - 1)^{3}\right|, 12(\operatorname{E}(\xi_{1}^{4}) - 1)\right)}{6} ||f^{(3)}||_{\infty} \sum_{l=1}^{n} \left(|\tilde{a}_{l,l}| \sum_{i=1}^{n} \tilde{a}_{i,l}^{2}\right)$$

$$+ \frac{16 \operatorname{E}(\xi_{1}^{8}) + 80 \operatorname{E}(\xi_{1}^{4}) + 3\tau^{4} + 96}{24} ||f^{(4)}||_{\infty} \sum_{l=1}^{n} \left(\sum_{i=1}^{n} \tilde{a}_{i,l}^{2}\right)^{2}.$$

This completes the proof.

Proof of Theorem 3. Throughout the proof, we use the similar notations as in Theorem 2 and define

$$S = \frac{(\sqrt{\phi} \boldsymbol{\epsilon})^{\top} \mathbf{A} \sqrt{\phi} \boldsymbol{\epsilon} - \text{tr}(\mathbf{A})}{\sqrt{2 \text{tr}(\mathbf{A}^2) + (\phi^2 \operatorname{E}(\epsilon_1^4) - 3) \text{tr}(\mathbf{A}^{\circ 2})}}$$

and

$$S_{\hat{\tau}}^* = \frac{\hat{\tau} \sum_{i=1}^n a_{i,i} \check{z}_i + 2 \sum_{1 \le i < j \le n} a_{i,j} z_i z_j}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + (\phi^2 \operatorname{E}(\epsilon_1^4) - 3) \operatorname{tr}(\mathbf{A}^{\circ 2})}},$$

where $z_1, \ldots, z_n, \check{z}_1, \ldots, \check{z}_n$ are iid $\mathcal{N}(0,1)$ random variables and are independent of $\hat{\tau}^2$.

By a standard subsequence argument, we only need to prove the theorem along a subsequence of $\{n\}$. Hence, without loss of generality, we assume $\hat{\tau}^2 \xrightarrow{a.s.} \phi^2 \operatorname{E}(\epsilon_1^4) - 1$. Write

$$\begin{split} S_{\hat{\tau}}^* = & \frac{\sqrt{\phi^2 \operatorname{E}(\epsilon_1^4) - 1} \sum_{i=1}^n a_{i,i} \check{z}_i + 2 \sum_{1 \leq i < j \leq n} a_{i,j} z_i z_j}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + (\phi^2 \operatorname{E}(\epsilon_1^4) - 3) \operatorname{tr}(\mathbf{A}^{\circ 2})}} \\ + & \frac{(\hat{\tau} - \sqrt{\phi^2 \operatorname{E}(\epsilon_1^4) - 1}) \sum_{i=1}^n a_{i,i} \check{z}_i}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + (\phi^2 \operatorname{E}(\epsilon_1^4) - 3) \operatorname{tr}(\mathbf{A}^{\circ 2})}} \\ =: & S_{\hat{\tau},1}^* + S_{\hat{\tau},2}^*. \end{split}$$

Note that $S_{\hat{\tau},1}^*$ is independent of $\hat{\tau}$. Since $\mathrm{E}(S_{\hat{\tau},1}^{*2})=1$, the distributions $\mathcal{L}(S_{\hat{\tau},1}^*)$ are tight as $n\to\infty$. Hence, without loss of generality, we assume $\mathcal{L}(S_{\hat{\tau},1}^*)$ weakly converges to a limit distribution with distribution function $F^{\dagger}(x)$. Let S^{\dagger} be a random variable with distribution function $F^{\dagger}(x)$. By some algebra (See, e.g., [7, Proposition A.1.(iii)]), it can be shown that $\mathrm{E}(S_{\hat{\tau},1}^{*4})$ is uniformly bounded. Then $\mathcal{L}(S_{\hat{\tau},1}^{*2})$ is uniformly integrable. Hence $\mathrm{E}(S^{\dagger 2})=1$ and $F^{\dagger}(x)$ can not concentrate on a single point. Consequently, $F^{\dagger}(x)$ is continuous and is strict increasing for $x\in\{x:0< F(x)<1\}$; see [32] as well the remark made by A. N. Kolmogorov on that paper.

The condition (7) implies that $\mathrm{E}[S_{\hat{\tau},2}^{*2}|\hat{\tau}] \to 0$ almost surely. Then almost surely, $\mathcal{L}(S_{\hat{\tau}}^*|\hat{\tau}) \leadsto \mathcal{L}(S^{\dagger})$. Consequently, for every $f \in \mathscr{C}^4(\mathbb{R})$, we have $|\mathrm{E}[f(S_{\hat{\tau}}^*)|\hat{\tau}] - \mathrm{E}[f(S^{\dagger})]| \to 0$ almost surely. On the other hand, Theorem 2 and the condition (7) imply $|\mathrm{E}[f(S)]| \to 0$ almost surely. Thus, $|\mathrm{E}[f(S)]| \to 0$. That is, $\mathcal{L}(S) \leadsto \mathcal{L}(S^{\dagger})$.

Note that

$$x^{(1)} = \frac{F^{-1}(1 - \alpha; \mathbf{A}, \hat{\tau}) - \operatorname{tr}(\mathbf{X}_b^* \mathbf{X}_b^{*\top})^{-1}}{n - q}.$$

We need to deal with $F^{-1}(1-\alpha; \mathbf{A}, \hat{\tau})$. Since $\mathcal{L}(S_{\hat{\tau}}^*|\hat{\tau}) \rightsquigarrow \mathcal{L}(S^{\dagger})$ almost surely, the fact

$$\Pr\left(S_{\hat{\tau}}^* > \frac{F^{-1}(1-\alpha; \mathbf{A}, \hat{\tau})}{\sqrt{2\operatorname{tr}(\mathbf{A}^2) + (\phi^2\operatorname{E}(\epsilon_1^4) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})}} \middle| \hat{\tau}\right) = \alpha$$

implies that almost surely,

$$\frac{F^{-1}(1-\alpha; \mathbf{A}, \hat{\tau})}{\sqrt{2\operatorname{tr}(\mathbf{A}^2) + (\phi^2 \operatorname{E}(\epsilon_1^4) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})}} \to F^{\dagger - 1}(1-\alpha). \tag{15}$$

We also need the fact that

$$(\sqrt{\phi}\boldsymbol{\epsilon})^{\top} \tilde{\mathbf{U}}_a \tilde{\mathbf{U}}_a^{\top} (\sqrt{\phi}\boldsymbol{\epsilon}) = (1 + o_p(1))(n - q), \tag{16}$$

which is a consequence of

$$\mathrm{E}\left((\sqrt{\phi}\boldsymbol{\epsilon})^{\top}\tilde{\mathbf{U}}_{a}\tilde{\mathbf{U}}_{a}^{\top}(\sqrt{\phi}\boldsymbol{\epsilon})\right) = n - q, \quad \mathrm{Var}\left((\sqrt{\phi}\boldsymbol{\epsilon})^{\top}\tilde{\mathbf{U}}_{a}\tilde{\mathbf{U}}_{a}^{\top}(\sqrt{\phi}\boldsymbol{\epsilon})\right) = O(n - q).$$

The fact $S \rightsquigarrow S^{\dagger}$, the equations (15), (16) and Slutsky's theorem lead to

$$\Pr\left(T > x^{(1)}\right)$$

$$= \Pr\left(T > \frac{F^{-1}(1 - \alpha; \mathbf{A}, \hat{\tau}) - \operatorname{tr}\left((\mathbf{X}_{b}^{*}\mathbf{X}_{b}^{*\top})^{-1}\right)}{n - q}\right)$$

$$= \Pr\left((\sqrt{\phi}\boldsymbol{\epsilon})^{\top}\mathbf{A}(\sqrt{\phi}\boldsymbol{\epsilon}) > \frac{(\sqrt{\phi}\boldsymbol{\epsilon})^{\top}\tilde{\mathbf{U}}_{a}\tilde{\mathbf{U}}_{a}^{\top}(\sqrt{\phi}\boldsymbol{\epsilon})}{n - q}F^{-1}(1 - \alpha; \mathbf{A}, \hat{\tau})\right)$$

$$= \Pr\left(S > \frac{(\sqrt{\phi}\boldsymbol{\epsilon})^{\top}\tilde{\mathbf{U}}_{a}\tilde{\mathbf{U}}_{a}^{\top}(\sqrt{\phi}\boldsymbol{\epsilon})}{n - q}\frac{F^{-1}(1 - \alpha; \mathbf{A}, \hat{\tau})}{\sqrt{2\operatorname{tr}(\mathbf{A}^{2}) + (\phi^{2}\operatorname{E}(\epsilon_{1}^{4}) - 3)\operatorname{tr}(\mathbf{A}^{\circ 2})}}\right)$$

$$= \Pr\left(S > (1 + o_{P}(1))F^{-1}(1 - \alpha)\right)$$

$$\to \alpha.$$

This proves the theorem.

Proof of Proposition 1. From [1, Theorem 2.1], one can obtain the explicit forms of $\operatorname{Var}\left(\tilde{\boldsymbol{\epsilon}}^{\top}\left(\tilde{\mathbf{P}}_{a}\right)\tilde{\boldsymbol{\epsilon}}\right)$ and $\operatorname{Var}\left(\sum_{i=1}^{n}\tilde{\epsilon}_{i}^{4}\right)$ which involves the traces of certain matrices. Using [19, Theorem 5.5.1], one can see that the eigenvalues of these matrices are all bounded. Hence it can be deduced that $\operatorname{Var}(\tilde{\boldsymbol{\epsilon}}^{\top}\tilde{\mathbf{P}}_{a}\tilde{\boldsymbol{\epsilon}}) = O(n)$ and $\operatorname{Var}\left(\sum_{i=1}^{n}\tilde{\epsilon}_{i}^{4}\right) = O(n)$. Thus,

$$\tilde{\boldsymbol{\epsilon}}^{\top} \tilde{\mathbf{P}}_{a} \tilde{\boldsymbol{\epsilon}} = (n - q) \phi^{-1} + O_{P}(\sqrt{n}),$$

$$\sum_{i=1}^{n} \tilde{\epsilon}_{i}^{4} = 3\phi^{-2} \operatorname{tr}(\tilde{\mathbf{P}}_{a}^{\circ 2}) + \left(\operatorname{E}(\epsilon_{1}^{4}) - 3\phi^{-2} \right) \operatorname{tr}\left(\tilde{\mathbf{P}}_{a}^{\circ 2}\right)^{2} + O_{P}(\sqrt{n}).$$

It follows that

$$\hat{\tau}^2 = \phi^2 \operatorname{E}(\epsilon_1^4) - 1 + O_P \left(\frac{\sqrt{n}}{\operatorname{tr} \left(\tilde{\mathbf{P}}_a^{\circ 2} \right)^2} \right).$$

Let $\delta_{i,j} = 1$ if i = j and 0 if $i \neq j$. We have

$$n = \sum_{i=1}^{n} \sum_{j=1}^{n} \delta_{i,j}^{4}$$

$$= \sum_{i=1}^{n} \sum_{j=1}^{n} \left((\tilde{\mathbf{P}}_{a})_{i,j} + (\mathbf{P}_{a})_{i,j} \right)^{4}$$

$$\leq 8 \sum_{i=1}^{n} \sum_{j=1}^{n} \left((\tilde{\mathbf{P}}_{a})_{i,j} \right)^{4} + 8 \sum_{i=1}^{n} \sum_{j=1}^{n} (\mathbf{P}_{a})_{i,j}^{4}$$

$$\leq 8 \sum_{i=1}^{n} \sum_{j=1}^{n} \left((\tilde{\mathbf{P}}_{a})_{i,j} \right)^{4} + 8 \sum_{i=1}^{n} \sum_{j=1}^{n} (\mathbf{P}_{a})_{i,j}^{2}$$

$$= 8 \operatorname{tr} \left(\tilde{\mathbf{P}}_{a}^{\circ 2} \right)^{2} + 8q.$$

Then

$$\frac{\sqrt{n}}{\operatorname{tr}\left(\tilde{\mathbf{P}}_{a}^{\circ 2}\right)^{2}} = O\left(\frac{1}{\sqrt{n}}\right).$$

This completes the proof.

Proof of Proposition 2. Without loss of generality, we assume **A** is a diagonal matrix and $|b_1| \geq \cdots \geq |b_n|$. By a standard subsequence argument, we only need to prove the result along a subsequence. Hence we can assume $\lim_{n\to\infty} \|b\|^2 / \operatorname{tr}(\mathbf{A}^2) = c \in [0, +\infty]$. If c = 0, Lyapunov central limit theorem implies that

$$\frac{Z^{\top} \mathbf{A} Z + b^{\top} Z - \operatorname{tr}(\mathbf{A})}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + \|\mathbf{b}\|^2}} = (1 + o_P(1)) \frac{Z^{\top} \mathbf{A} Z - \operatorname{tr}(\mathbf{A})}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2)}} + o_P(1) \rightsquigarrow \mathcal{N}(0, 1).$$

If $c = +\infty$,

$$\frac{Z^{\top} \mathbf{A} Z + b^{\top} Z - \operatorname{tr}(\mathbf{A})}{\sqrt{2 \operatorname{tr}(\mathbf{A}^2) + \|\mathbf{b}\|^2}} = (1 + o_P(1)) \frac{b^{\top} Z}{\|\mathbf{b}\|} + o_P(1) \rightsquigarrow \mathcal{N}(0, 1).$$

In what follows, we assume $c \in (0, +\infty)$. By Helly selection theorem, we can assume $\lim_{n\to\infty} |b_i|/\|\mathbf{b}\| = b_i^* \in [0, 1], i = 1, 2, \ldots$ From Fatou's lemma, we have $\sum_{i=1}^{\infty} (b_i^*)^2 \leq 1$. Consequently, $\lim_{i\to\infty} b_i^* = 0$.

Note that the condition $\operatorname{tr}(\mathbf{A}^4)/\operatorname{tr}^2(\mathbf{A}^2) \to 0$ is equivalent to $\lambda_1(\mathbf{A}^2)/\operatorname{tr}(\mathbf{A}^2) \to 0$. Then for every fixed integer r > 0,

$$\frac{\sum_{i=1}^{r} a_{i,i}^2}{\sum_{i=1}^{n} a_{i,i}^2} \le \frac{r \max_{1 \le i \le n} a_{i,i}^2}{\sum_{i=1}^{n} a_{i,i}^2} \to 0.$$

Then there exists a sequence of positive integers $r(n) \to \infty$ such that $\left(\sum_{i=1}^{r(n)} a_{i,i}^2\right) / \left(\sum_{i=1}^n a_{i,i}^2\right) \to 0$ and $r(n)/n \to 0$. Write

$$Z^{\top} \mathbf{A} Z + b^{\top} Z - \operatorname{tr}(\mathbf{A}) = \sum_{i=1}^{r(n)} a_{i,i} (z_i^2 - 1) + \sum_{i=1}^{r(n)} b_i z_i + \sum_{i=r(n)+1}^{n} \left(a_{i,i} (z_i^2 - 1) + b_i z_i \right),$$

which is a sum of independent random variables. The first term is negligible since $\operatorname{Var}(\sum_{i=1}^{r(n)} a_{i,i}(z_i^2 - 1)) = o(\sum_{i=1}^{n} a_{i,i}^2)$. Now we deal with the third term. From Berry-Esseen inequality (See, e.g., [10, Theorem 11.2]), there exists an absolute constant $C^* > 0$, such that

$$\sup_{x \in \mathbb{R}} \left| \Pr\left(\frac{\sum_{i=r(n)+1}^{n} \left(a_{i,i}(z_{i}^{2}-1) + b_{i}z_{i} \right)}{\sqrt{2 \sum_{i=r(n)+1}^{n} a_{i,i}^{2} + \sum_{i=r(n)+1}^{n} b_{i}^{2}}} \le x \right) - \Phi(x) \right|$$

$$\le C^{*} \frac{\sum_{i=r(n)+1}^{n} \operatorname{E} \left| a_{i,i}(z_{i}^{2}-1) + b_{i}z_{i} \right|^{3}}{\left(2 \sum_{i=r(n)+1}^{n} a_{i,i}^{2} + \sum_{i=r(n)+1}^{n} b_{i}^{2} \right)^{3/2}}.$$

By some simple algebra, there exist absolute constants $C_1^*, C_2^* > 0$ such that for sufficiently large n,

$$\sup_{x \in \mathbb{R}} \left| \Pr\left(\frac{\sum_{i=r(n)+1}^{n} \left(a_{i,i}(z_{i}^{2}-1) + b_{i}z_{i} \right)}{\sqrt{2 \sum_{i=r(n)+1}^{n} a_{i,i}^{2} + \sum_{i=r(n)+1}^{n} b_{i}^{2}}} \le x \right) - \Phi(x) \right|$$

$$\le C_{1}^{*} \frac{\max_{1 \le i \le n} |a_{i,i}|}{\sqrt{\sum_{i=1}^{n} a_{i,i}^{2}}} + C_{2}^{*} \frac{|b_{r(n)+1}|}{\|\mathbf{b}\|}.$$

Since the right hand side tends to 0, we have

$$\frac{\sum_{i=r(n)+1}^{n} \left(a_{i,i}(z_i^2 - 1) + b_i z_i \right)}{\sqrt{2 \sum_{i=r(n)+1}^{n} a_{i,i}^2 + \sum_{i=r(n)+1}^{n} b_i^2}} \rightsquigarrow \mathcal{N}(0,1).$$

Note that $\sum_{i=r(n)+1}^{n} (a_{i,i}(z_i^2 - 1) + b_i z_i)$ is independent of $\sum_{i=1}^{r(n)} b_i z_i$ and $\sum_{i=1}^{r(n)} b_i z_i \sim \mathcal{N}(0, \sum_{i=1}^{r(n)} b_i^2)$. Thus,

$$\frac{\sum_{i=1}^{r(n)} b_i z_i + \sum_{i=r(n)+1}^n \left(a_{i,i} (z_i^2 - 1) + b_i z_i \right)}{\sqrt{2 \sum_{i=1}^n a_{i,i}^2 + \sum_{i=1}^n b_i^2}} \rightsquigarrow \mathcal{N}(0, 1).$$

This completes the proof.

Note that under the normality, $T_n - \operatorname{tr}((\mathbf{X}_b^* \mathbf{X}_b^{*\top})^{-1})/(n-q)$ has zero mean.

Proof of Theorem 4. We note that

$$\Pr\left(\frac{\mathbf{y}^{*\top} \left(\mathbf{X}_{b}^{*} \mathbf{X}_{b}^{*\top}\right)^{k} \mathbf{y}^{*}}{\mathbf{y}^{*\top} \mathbf{y}^{*}} \leq \mathrm{E}(\gamma_{I}^{k}) + \sqrt{\frac{2 \operatorname{Var}\left(\gamma_{I}^{k}\right)}{n - q}} x\right)$$

$$= \Pr\left(\mathbf{y}^{*\top} \left(\mathbf{X}_{b}^{*} \mathbf{X}_{b}^{*\top}\right)^{k} \mathbf{y}^{*} \leq \left(\mathrm{E}(\gamma_{I}^{k}) + \sqrt{\frac{2 \operatorname{Var}\left(\gamma_{I}^{k}\right)}{n - q}} x\right) \mathbf{y}^{*\top} \mathbf{y}^{*}\right)$$

$$= \Pr\left(\mathbf{y}^{*\top} \mathbf{B} \mathbf{y}^{*} \leq 0\right),$$
(17)

where

$$\mathbf{B} = \left(\mathbf{X}_b^* \mathbf{X}_b^{*\top}\right)^k - \left(\mathbf{E}(\gamma_I^k) + \sqrt{\frac{2 \operatorname{Var}\left(\gamma_I^k\right)}{n - q}} x\right) \mathbf{I}_{n - q}.$$

Since $\mathbf{y}^{*\top}\mathbf{B}\mathbf{y}^{*} = \boldsymbol{\epsilon}^{\top}\tilde{\mathbf{U}}_{a}\mathbf{B}\tilde{\mathbf{U}}_{a}^{\top}\boldsymbol{\epsilon} + 2\boldsymbol{\epsilon}^{\top}\tilde{\mathbf{U}}_{a}\mathbf{B}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b} + \boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b}$, we have

$$\Pr\left(\mathbf{y}^{*\top}\mathbf{B}\mathbf{y}^{*} \leq 0\right)$$

$$=\Pr\left(\frac{\boldsymbol{\epsilon}^{\top}\tilde{\mathbf{U}}_{a}\mathbf{B}\tilde{\mathbf{U}}_{a}^{\top}\boldsymbol{\epsilon} + 2\boldsymbol{\epsilon}^{\top}\tilde{\mathbf{U}}_{a}\mathbf{B}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b} - \phi^{-1}\operatorname{tr}(\mathbf{B})}{\sqrt{2\phi^{-2}\operatorname{tr}(\mathbf{B}^{2}) + 4\phi^{-1}\boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}^{2}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b}}}\right)$$

$$\leq \frac{-\boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b} - \phi^{-1}\operatorname{tr}(\mathbf{B})}{\sqrt{2\phi^{-2}\operatorname{tr}(\mathbf{B}^{2}) + 4\phi^{-1}\boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}^{2}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b}}}\right).$$

To apply proposition 2, we need to verify the condition $\lambda_1\left(\mathbf{B}^2\right)/\operatorname{tr}\left(\mathbf{B}^2\right)\to 0$. It is straightforward to show that $\operatorname{tr}(\mathbf{B}^2)=(n-q+2x^2)\operatorname{Var}(\gamma_I^k)$. On the other hand,

$$\lambda_1 \left(\mathbf{B}^2 \right) = \max_{1 \le i \le n - q} \left(\gamma_i^k - \mathrm{E}(\gamma_I^k) - \sqrt{\frac{2 \operatorname{Var} \left(\gamma_I^k \right)}{n - q}} x \right)^2$$
$$\le 2 \max_{1 \le i \le n - q} \left(\gamma_i^k - \mathrm{E}(\gamma_I^k) \right)^2 + 4 \frac{\operatorname{Var} \left(\gamma_I^k \right)}{n - q} x^2.$$

Thus,

$$\begin{split} \frac{\lambda_{1}(\mathbf{B}^{2})}{\operatorname{tr}(\mathbf{B}^{2})} \leq & 2\frac{\max_{1 \leq i \leq n-q} \left(\gamma_{i}^{k} - \operatorname{E}(\gamma_{I}^{k})\right)^{2}}{(n-q+2x^{2})\operatorname{Var}(\gamma_{I}^{k})} + 4\frac{x^{2}}{(n-q)(n-q+x^{2})} \\ \leq & 2\frac{\max_{1 \leq i \leq n-q} \left(\gamma_{i}^{k} - \operatorname{E}(\gamma_{I}^{k})\right)^{2}}{(n-q)\operatorname{Var}(\gamma_{I}^{k})} + \frac{4}{(n-q)}, \end{split}$$

which tends to 0 by the condition (8). Hence Proposition 2 implies that

$$\Pr\left(\mathbf{y}^{*\top}\mathbf{B}\mathbf{y}^{*} \leq 0\right) = \Phi\left(\frac{-\boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b} - \phi^{-1}\operatorname{tr}(\mathbf{B})}{\sqrt{2\phi^{-2}\operatorname{tr}(\mathbf{B}^{2}) + 4\phi^{-1}\boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}^{2}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b}}}\right) + o(1).$$
(18)

Then the conclusion follows from (17), (18) and the following facts

$$\operatorname{tr}(\mathbf{B}) = -(n-q)\sqrt{\frac{2\operatorname{Var}\left(\gamma_{I}^{k}\right)}{n-q}}x,$$

$$\operatorname{tr}(\mathbf{B}^{2}) = (1+o(1))(n-q)\operatorname{Var}(\gamma_{I}^{k}),$$

$$\boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b} = (n-q)\left(\operatorname{Cov}\left(\gamma_{I}^{k},\gamma_{I}w_{I}^{2}\right) - \operatorname{E}(\gamma_{I}w_{I}^{2})\sqrt{\frac{2\operatorname{Var}\left(\gamma_{I}^{k}\right)}{n-q}}x\right),$$

$$\boldsymbol{\beta}_{b}^{\top}\mathbf{X}_{b}^{*\top}\mathbf{B}^{2}\mathbf{X}_{b}^{*}\boldsymbol{\beta}_{b} = (n-q)\operatorname{E}\left[\left(\gamma_{I}^{k} - \operatorname{E}(\gamma_{I}^{k}) - \sqrt{\frac{2\operatorname{Var}(\gamma_{I}^{k})}{n-q}}x\right)^{2}\gamma_{I}w_{I}^{2}\right].$$

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